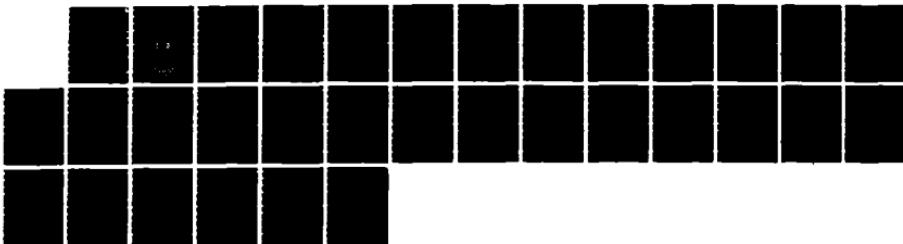
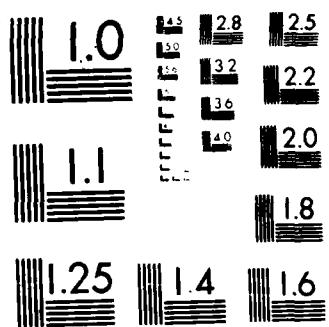


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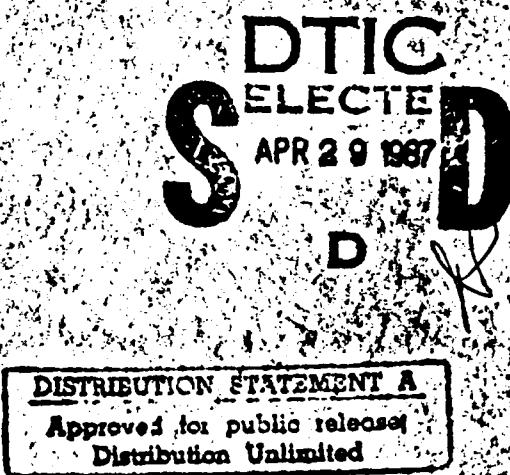
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Recursive Linear Smoothing
for the 2-D Helmholtz Equation

Laurence R. Riddle and Howard L. Weinst

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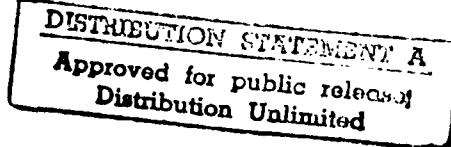
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Recursive Linear Smoothing for the 2-D Helmholtz Equation

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ABSTRACT

A fast algorithm for reconstructing images governed by a 2-D Helmholtz equation is presented. The computational complexity of the algorithm is $O(NM \log M)$ or $O(NM^2)$ depending on boundary conditions, where N and M are the number of spatial grid points in the x and y directions respectively. This problem arises when smoothing a large number of images governed by the 2-D wave equation, because a Fourier transform in time gives a new set of images governed by the Helmholtz equation. When the images come from a scattering process, we show that a linear least-squares Born inversion of the wave field amplitudes can be performed during the smoothing procedure without changing the computational complexity. We also show that the smoothing algorithm is well-posed, and that the sample functions of the smoothed estimate possess smoothness properties consistent with the Helmholtz equation.

1. Introduction

In this paper we derive a fast, recursive, linear least-squares smoothing algorithm for the 2-D Helmholtz equation. Our algorithm can be used, for example, to smooth a large number of images governed by the 2-D wave equation arising in acoustical holography [8] or in oceanic surveillance. If we assume that the vibrating system is in steady state, and that the inputs and observation noise are temporally stationary, then a

* This work was supported by the Office of Naval Research under Contract N00014-85-K-0255



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Fourier expansion with respect to time gives a new set of images, indexed by the frequency f , that are governed by the Helmholtz equation whose wavenumber k varies with f . Each transformed image can be smoothed separately, and estimates of the original time-domain images can then be obtained by an inverse Fourier transform.

In order to smooth images governed by the Helmholtz equation, we reformulate the equation as a well-posed, distributed-parameter, acausal linear system, and use the recent extension of Adams, Willsky and Levy [1] of the method of complementary models [9] to write the Hamiltonian system for the smoothed estimate. Transforming in one direction produces a set of indexed, well-posed, finite-dimensional, acausal linear systems which we solve recursively using a diagonalizing change of variables. The complexity of our algorithm for each wavenumber k is $O(NM\log M)$ or $O(NM^2)$ depending on the boundary conditions, where N is the number of grid points in the x direction and M is the number in the (orthogonal) y direction.

Another approach to the smoothing of images governed by the Helmholtz equation involves the use of the Karhunen-Loeve expansion of the wave field [2],[3],[4], resulting in an algorithm with complexity $O(MN\log MN)$. However, the boundary conditions are required to be conservative and have no random inputs.

In contrast, we can handle x-boundary conditions that are conservative or dissipative and that include random inputs. We note that work using a similar approach is discussed by Yoshida and Ogura in [10]. In their work, the dynamics are discrete and the underlying random field is homogeneous (spatially stationary), whereas in this *paper*, the dynamics are continuous and the random field is not required to be homogeneous. Furthermore, an important step in the derivation of the estimator in [10] is the replacement of a non-Markovian random process with a Markovian random process having the same mean and covariance. In our approach, this realization step is not needed.

2. Problem Statement

Consider the scalar Helmholtz equation on the rectangle $\Omega = [0, L_1] \times [0, L_2]$:

$$u_{xx} + u_{yy} + k^2 u = \epsilon(x, y) \quad (2.1)$$

$$u = u(x, y), (x, y) \in \Omega$$

with boundary conditions

$$u(0, y) = v_1(y), \quad u(L_1, y) = v_2(y) \quad (2.2a)$$

$$u(x, 0) = 0, \quad u(x, L_2) = 0 \quad (2.2b)$$

Here

$$k^2 = k_0^2 + j\eta, \quad k_0 = \frac{2\pi f}{c}, \quad \eta > 0$$

ϵ is the input field, v_1 and v_2 are boundary inputs on the x-axis, u is the wavefield amplitude, f is the temporal

frequency of interest, c is the phase velocity of the medium, and η is the damping term. The observations are

$$z(x,y) = u(x,y) + w(x,y), (x,y) \in \Omega \quad (2.3)$$

where $w(x,y)$ is the observation noise field.

We shall make the following statistical assumptions:

(1) the driving field ϵ and observation noise field w are zero mean and white with constant intensities q and r , respectively, and are uncorrelated with each other, (2) if $v(y) = [v_1(y) \ v_2(y)]'$ then

$$Ev(y) = 0$$

$$Ev(y)v(s)' = \Pi_v \delta(y-s)$$

where Π_v is invertible and v is uncorrelated with ϵ and w .

The estimation problem of interest here is to determine the linear least-squares estimate $\hat{u}(x,y)$ of $u(x,y)$, $(x,y) \in \Omega$, given the observations (2.3) over the entire rectangle Ω .

To see how the Helmholtz equation may arise in practice, consider the 2-D wave equation on the rectangle $\Omega = [0, L_1] \times [0, L_2]$:

$$u_{tt} - c^2(u_{xx} + u_{yy}) + \gamma u_t = d(x,y,t)$$

$$u = u(x,y,t), (x,y,t) \in \Omega \times [T_0, T_1]$$

with boundary conditions

$$u(0,y,t) = v_1(y,t), \quad u(L_1,y,t) = v_2(y,t)$$

$$u(x,0,t) = 0, \quad u(x,L_2,t) = 0$$

$$u(x,y,T_0) = u_t(x,y,T_0) = 0$$

and observations

$$z(x,y,t) = u(x,y,t) + w(x,y,t)$$

Let $T_0 \rightarrow -\infty$ and assume that the observation interval is $[0, T_1]$ where T_1 is very large. Then a Fourier series expansion of the observations will give a new set of images $z(x,y,f)$ that are governed by (2.1)-(2.3), where the dependence on f has been suppressed, and where $\epsilon = -d/c^2$ and $\eta = 2\pi f \gamma$. We assume that the input field $d(x,y,t)$, boundary inputs $v_1(y,t)$, $v_2(y,t)$, and observation noise $w(x,y,t)$ are wide-sense stationary in time. The stationarity assumption implies that for $f_1 \neq f_2$, $z(x,y,f_1)$ and $u(x,y,f_1)$ are uncorrelated with $z(x,y,f_2)$, since $T_1 \rightarrow \infty$. We can therefore solve an uncoupled set of smoothing problems for the 2-D Helmholtz equation (indexed by f), and then inverse transform to recover the time-domain estimates.

3. State Space Formulation and Characterization of the Estimate

In order to put (2.1)-(2.3) in state-variable form, define an operator T with domain $D(T)$ as follows:

$$Tf = -(k^2 f(x,y) + f_{yy}(x,y)), \quad f \in D(T)$$

where

$$D(T) = \{f \in L_2(\Omega) : f, f_y \text{ abs. cont.}, f_{yy} \in L_2(\Omega), \\ f(x,0) = f(x,L_2) = 0\}$$

Also define the state vector $m(x,y)$ as

$$m(x,y) = [m_1(x,y), m_2(x,y)]' = [u(x,y), u_x(x,y)]'$$

We can now rewrite (2.1)-(2.3) as

$$\frac{\partial}{\partial x} m(x,y) = Am(x,y) + B\epsilon(x,y) \quad (3.1a)$$

$$m_1 \in D(T)$$

$$v(y) = V_0 m(0,y) + V_L m(L_1,y) \quad (3.1b)$$

$$z(x,y) = Cm(x,y) + w(x,y) \quad (3.1c)$$

where

$$A = \begin{bmatrix} 0 & 1 \\ T & 0 \end{bmatrix}, \quad B = \begin{bmatrix} 0 \\ 1 \end{bmatrix}, \quad C = \begin{bmatrix} 1 & 0 \end{bmatrix}, \quad v = \begin{bmatrix} v_1 \\ v_2 \end{bmatrix}$$

$$V_0 = \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix}, \quad V_L = \begin{bmatrix} 0 & 0 \\ 1 & 0 \end{bmatrix}$$

Equations (3.1) are in the form of a distributed parameter, acausal linear system. Finite-dimensional acausal linear systems are discussed in [6]-[7]. We show in Appendix A that (3.1) are well posed, in the sense that if $\epsilon = v = 0$, then $m = 0$.

Using results in [1], we can express the linear least-squares estimate \hat{m} of m , given the observations (2.3), as the solution of the following Hamiltonian system:

$$\frac{\partial}{\partial x} \begin{bmatrix} \hat{m}(x,y) \\ \hat{\lambda}(x,y) \end{bmatrix} = \begin{bmatrix} A & qBB' \\ r^{-1}C'C - A^* & \end{bmatrix} \begin{bmatrix} \hat{m}(x,y) \\ \hat{\lambda}(x,y) \end{bmatrix} + \begin{bmatrix} 0 \\ C'r^{-1}z(x,y) \end{bmatrix} \quad (3.2a)$$

$$\hat{m} \in D(A), \hat{\lambda} \in D(A^*)$$

$$0 = V^* \Pi_v^{-1} V \begin{bmatrix} \hat{m}(0, y) \\ \hat{m}(L_1, y) \end{bmatrix} + \begin{bmatrix} -I & 0 \\ 0 & I \end{bmatrix} \begin{bmatrix} \hat{\lambda}(0, y) \\ \hat{\lambda}(L_1, y) \end{bmatrix} \quad (3.2b)$$

where

$$V = [V_0 \ V_L]$$

the asterisk denotes adjoint, and $D(A), D(A^*)$ denote the domains of A and A^* . These domains are

$$D(A) = \{[f_1, f_2]': f_1 \in D(T), f_2 \in \mathbf{L}_2(\Omega)\}$$

$$D(A^*) = \{[f_1, f_2]': f_1 \in \mathbf{L}_2(\Omega), f_2 \in D(T)\}$$

Moreover, the input estimate satisfies

$$\hat{\epsilon}(x, y) = qB^* \hat{\lambda}(x, y) \quad (3.2c)$$

Eq. (3.2c) can be interpreted as a linear least-squares Born inversion when the observations are the scattered wavefield in an inverse scattering experiment.

Instead of solving (3.2) directly for \hat{m} and $\hat{\lambda}$, we will transform (3.2) with respect to y using the discrete sine transform \mathbf{S} , given by

$$\mathbf{S}g = \frac{1}{L_2} \int_0^{L_2} \sin(py) g(y) dy$$

$$p = \frac{2\pi n}{L_2}, n = 0, \pm 1, \pm 2, \dots$$

It can easily be verified that $\mathbf{S}T = (p^2 - k^2)\mathbf{S}$ and thus $\mathbf{S}A = A_p \mathbf{S}$ where

$$A_p = \begin{bmatrix} 0 & 1 \\ p^2 - k^2 & 0 \end{bmatrix}$$

Transforming (3.2a) with respect to y gives

$$\frac{\partial}{\partial x} \begin{bmatrix} \hat{m}(x, p) \\ \hat{\lambda}(x, p) \end{bmatrix} = \begin{bmatrix} A_p & qBB' \\ r^{-1}C'C & -A_p^* \end{bmatrix} \begin{bmatrix} \hat{m}(x, p) \\ \hat{\lambda}(x, p) \end{bmatrix} + \begin{bmatrix} 0 \\ C'r^{-1}z(x, p) \end{bmatrix} \quad (3.3a)$$

where $\hat{m}(x, p) = S\hat{m}(x, y)$, etc. Transforming (3.2b) with respect to y gives

$$0 = V^* \Pi_v^{-1} V \begin{bmatrix} \hat{m}(0, p) \\ \hat{m}(L_1, p) \end{bmatrix} + \begin{bmatrix} -I & 0 \\ 0 & I \end{bmatrix} \begin{bmatrix} \hat{\lambda}(0, p) \\ \hat{\lambda}(L_1, p) \end{bmatrix}.$$

When written in the standard form for an acausal linear system, these boundary conditions are

$$0 = W_0 \begin{bmatrix} \hat{m}(0, p) \\ \hat{\lambda}(0, p) \end{bmatrix} + W_L \begin{bmatrix} \hat{m}(L_1, p) \\ \hat{\lambda}(L_1, p) \end{bmatrix} \quad (3.3b)$$

where

$$W_0 = \begin{bmatrix} V_0^* \Pi_v^{-1} V_0 & -I \\ V_L^* \Pi_v^{-1} V_0 & 0 \end{bmatrix} \text{ and } W_L = \begin{bmatrix} V_0^* \Pi_v^{-1} V_L & 0 \\ V_L^* \Pi_v^{-1} V_L & I \end{bmatrix}$$

We see that in the p domain the estimate Hamiltonian is decoupled; that is, (3.3) are simply indexed by p , and can be solved individually. In the next section we show how to solve (3.3) recursively for each p . Before doing so, we first consider whether (3.3) is well-posed. These equations can be shown to be well-posed by first realizing that for each p , (3.3) is the estimate Hamiltonian of another acausal linear system, the so-called p -dynamics of the Helmholtz equation.

$$\frac{\partial}{\partial x} m(x, p) = A_p m(x, p) + B \epsilon(x, p) \quad (3.4a)$$

$$V_0 m(0, p) + V_L m(L_1, p) = v(p) \quad (3.4b)$$

$$z(x,p) = Cm(x,p) + w(x,p) \quad (3.4c)$$

where we have transformed (3.1) into the p -domain by S . The invertibility of $V_0 + V_L e^{A_p L_1}$ is necessary and sufficient for (3.4) to be well-posed [6]-[7]. It is easy to verify that this matrix is invertible for all p . To prove then that (3.3) is well-posed for all p , assume that the input $z(x,p)$ is identically zero. Now $\hat{m}(x,p)$ is the linear least-squares estimate of $m(x,p)$ based on observations which are identically zero, so $\hat{m}(x,p) = Em(x,p) = 0$, the last equality following from the well-posedness of (3.1).

As a result $\hat{\lambda}(x,p)$ satisfies

$$\frac{\partial}{\partial x} \hat{\lambda}(x,p) = -A_p^* \hat{\lambda}(x,p)$$

and

$$\hat{\lambda}(0,p) = \hat{\lambda}(L_1, p) = 0$$

and thus $\hat{\lambda}(x,p) = 0$.

4. Recursive Solution of the Estimate Hamiltonian

We shall derive our recursive smoothing formula by diagonalizing the dynamics in (3.3a) via a change of variables. Let

$$H_p = \begin{bmatrix} A_p & qBB' \\ r^{-1}C'C & -A_p^* \end{bmatrix}$$

The characteristic equation for H_p is given by

$$\lambda^4 + 2\operatorname{Re}(\alpha)\lambda^2 + |\alpha|^2 + r^{-1}q = 0$$

where $\alpha = p^2 - k^2$. Solutions to this equation are $\lambda_0, \overline{\lambda_0}, -\lambda_0, -\overline{\lambda_0}$, where

$$\lambda_0 = \left(\frac{(|\alpha|^2 + r^{-1}q)^{1/2} + \operatorname{Re}\alpha}{2} \right)^{1/2} - j \left(\frac{(|\alpha|^2 + r^{-1}q)^{1/2} - \operatorname{Re}\alpha}{2} \right)^{1/2} \quad (4.1)$$

The four eigenvalues of H_p are all distinct because both the real and imaginary parts of (4.1) are non-zero for all p . As a result, we can diagonalize H_p as follows :

$$M_p^{-1} H_p M_p = \begin{bmatrix} \Lambda_p & 0 \\ 0 & -\Lambda_p \end{bmatrix}$$

where

$$\Lambda_p = \begin{bmatrix} \lambda_0 & 0 \\ 0 & \overline{\lambda_0} \end{bmatrix}$$

and the modal matrix M_p and its inverse are given explicitly by

$$M_p = \begin{bmatrix} 1 & 1 & 1 & 1 \\ \lambda_0 & \overline{\lambda_0} & -\lambda_0 & -\overline{\lambda_0} \\ -\lambda_0 c & -\overline{\lambda_0} d & \lambda_0 c & \overline{\lambda_0} d \\ c & d & c & d \end{bmatrix}$$

$$M_p^{-1} = \frac{jq}{4\sigma} \begin{bmatrix} d & d\lambda_0^{-1} & \lambda_0^{-1} & -1 \\ -c & -c\overline{\lambda_0}^{-1} & -\overline{\lambda_0}^{-1} & 1 \\ d & -d\lambda_0^{-1} & -\lambda_0^{-1} & -1 \\ -c & c\overline{\lambda_0}^{-1} & \overline{\lambda_0}^{-1} & 1 \end{bmatrix}$$

with

$$c = j(\eta + \sigma)/q$$

$$d = j(\eta - \sigma)/q$$

$$\sigma = (qr^{-1} + \eta^2)^{1/2}$$

Now by using the following change of variables in (3.3):

$$\begin{bmatrix} \Psi_f(x, p) \\ \Psi_b(x, p) \end{bmatrix} = M_p^{-1} \begin{bmatrix} \hat{m}(x, p) \\ \hat{\lambda}(x, p) \end{bmatrix} \quad (4.2)$$

we get

$$\frac{\partial}{\partial x} \begin{bmatrix} \Psi_f \\ \Psi_b \end{bmatrix} = \begin{bmatrix} \Lambda_p & 0 \\ 0 & -\Lambda_p \end{bmatrix} \begin{bmatrix} \Psi_f \\ \Psi_b \end{bmatrix} + \begin{bmatrix} B_f \\ B_b \end{bmatrix} z \quad (4.3a)$$

$$0 = \begin{bmatrix} V_f^0 & V_b^0 \end{bmatrix} \begin{bmatrix} \Psi_f(0, p) \\ \Psi_b(0, p) \end{bmatrix} + \begin{bmatrix} V_f^L & V_b^L \end{bmatrix} \begin{bmatrix} \Psi_f(L_1, p) \\ \Psi_b(L_1, p) \end{bmatrix} \quad (4.3b)$$

where

$$\begin{bmatrix} V_f^0 & V_b^0 \end{bmatrix} = W_0 M_p, \quad \begin{bmatrix} V_f^L & V_b^L \end{bmatrix} = W_L M_p$$

$$\begin{bmatrix} B_f \\ B_b \end{bmatrix} = M_p^{-1} \begin{bmatrix} 0 \\ C'r^{-1} \end{bmatrix}$$

Equations (4.3) are in a form which can be solved recursively. In terms of $\Psi_f(0, p)$ and $\Psi_b(L_1, p)$, a solution to (4.3a) is

$$\Psi_f(x, p) = e^{\Lambda_p(x)} \Psi_f(0, p) + \Psi_f^0(x, p) \quad (4.4a)$$

$$\Psi_b(x, p) = e^{\Lambda_p(L_1-x)} \Psi_b(L_1, p) + \Psi_b^0(x, p) \quad (4.4b)$$

where

$$\frac{\partial}{\partial x} \Psi_f^0(x, p) = \Lambda_p \Psi_f^0(x, p) + B_f z(x, p) \quad (4.5a)$$

$$\frac{\partial}{\partial x} \Psi_b^0(x, p) = -\Lambda_p \Psi_b^0(x, p) + B_b z(x, p) \quad (4.5b)$$

$$\Psi_f^0(0, p) = 0, \quad \Psi_b^0(L_1, p) = 0. \quad (4.5c)$$

Note that (4.5a) and (4.5b) are stable in the forward and backward directions, respectively. Setting $x = L_1$ in

(4.4a), and $x = 0$ in (4.4b), we can solve for $\Psi_f(0, p)$ and $\Psi_b(L_1, p)$ in (4.3b) as

$$\begin{bmatrix} \Psi_f(0, p) \\ \Psi_b(L_1, p) \end{bmatrix} = -F_{fb}^{-1}\{V_f^L\Psi_f^0(L_1, p) + V_b^0\Psi_b^0(0, p)\}$$

where

$$F_{fb} = [V_f^0 + V_f^L e^{\Lambda, L_1} : V_b^L + V_b^0 e^{\Lambda, L_1}]$$

A solution to (4.3) is therefore given by

$$\begin{bmatrix} \Psi_f(x, p) \\ \Psi_b(x, p) \end{bmatrix} = \begin{bmatrix} -e^{\Lambda, x} & 0 \\ 0 & -e^{\Lambda, (L_1-x)} \end{bmatrix} F_{fb}^{-1}\{V_f^L\Psi_f^0(L_1, p) + V_b^0\Psi_b^0(0, p)\} \\ + \begin{bmatrix} \Psi_f^0(x, p) \\ \Psi_b^0(x, p) \end{bmatrix} \quad (4.6)$$

The F_{fb} matrix is invertible because

$$F_{fb} = F_H M_p \begin{bmatrix} I & 0 \\ 0 & e^{\Lambda, L_1} \end{bmatrix},$$

where F_H is the invertible matrix associated with (3.3) being well-posed:

$$F_H = W_0 + W_L e^{H, L_1}.$$

The behavior of the algorithm as $p \rightarrow \infty$ needs to be examined further. We will show that the determinant of F_{fb} does not vanish as $p \rightarrow \infty$. As p gets large, one can ignore the exponential terms in F_{fb} and write

$$F_{fb} \approx \begin{bmatrix} \theta_{11} + \lambda_0 c & \theta_{11} + \bar{\lambda}_0 d & \theta_{12} & \theta_{12} \\ -c & -d & 0 & 0 \\ \theta_{21} & \theta_{21} & \theta_{22} + \lambda_0 c & \theta_{22} + \bar{\lambda}_0 d \\ 0 & 0 & c & d \end{bmatrix}$$

where θ_{ij} is the ij -th entry of Π_v^{-1} . Then

$$\det F_{fb} \approx -(c-d)^2 \det(\Pi_v^{-1}) - c^2 d^2 (\lambda_0 - \bar{\lambda}_0)^2 + cd(c-d)(\lambda_0 - \bar{\lambda}_0)(\theta_{11} + \theta_{22})$$

As $p \rightarrow \infty$, $\lambda_0 - \bar{\lambda}_0 \rightarrow 0$, so that

$$\lim_{p \rightarrow \infty} \det F_{fb} = 4(qr^{-1} + \eta^2)q^{-2} \det(\Pi_v^{-1}) \neq 0$$

It is shown in Appendix B that under realistic energy assumptions on the observed images, $\Psi_f(x,p)$ and $\Psi_b(x,p)$ decay fast enough as $p \rightarrow \infty$ to ensure that $\hat{m}(x,y) \in D(A)$ and $\hat{\lambda} \in D(A^*)$.

To summarize, the solution procedure for solving the estimator equations (3.2) is (1) transform the observations into the p -domain, (2) compute $\Psi_f^0(x,p)$ and $\Psi_b^0(x,p)$ using (4.5), (3) find $\Psi_f(x,p)$ and $\Psi_b(x,p)$ from (4.6), (4) compute $\hat{m}(x,p)$ using (4.2), (5) inverse transform $\hat{m}(x,p)$ to get $\hat{m}(x,y)$.

5. Other Boundary Conditions

The smoothing problem for the 2-D Helmholtz equation with more general boundary conditions can be handled in a manner very similar to that just described. In addition, a wavenumber that varies in the y direction can be studied. The boundary conditions we consider in this section are conservative on the y -boundaries, and, in general, dissipative on the x -boundaries. To derive the estimator for such problems, we start off with the same 2-D Helmholtz equation (2.1), and x -boundary conditions (2.2a). The y -boundary conditions are now

$$\cos \beta u(x,0) + \sin \beta u_y(x,0) = 0 \quad (5.1a)$$

$$\cos\gamma u(x, L_2) + \sin\gamma u_y(x, L_2) = 0 \quad (5.1b)$$

where β and γ are real. Periodic boundary conditions

$$u(x, 0) = u(x, L_2), \quad u_y(x, 0) = u_y(x, L_2) \quad (5.2)$$

could also be assumed. The wavenumber k satisfies

$$k^2(y) = k_0^2(y) + j\eta, \quad \eta > 0$$

The estimate Hamiltonian for these boundary conditions and wavenumber is the same as (3.2), however the operator T is defined as

$$Tf = -(\partial_{yy} + k^2(y))f$$

with domain

$$D(T) = \{f \in L_2(\Omega) : f, f_y \text{ abs. cont., } f_{yy} \in L_2(\Omega), \\ f \text{ satisfies (5.1) or (5.2)}\}$$

We now introduce the selfadjoint Sturm-Liouville operator

$$Qf = -(\partial_{yy} + k_0^2(y))f$$

with domain

$$D(Q) = D(T)$$

The operator Q has a countably infinite number of eigenvalues μ_p and eigenfunctions $\Phi_p(y)$. If we use the transform operator \mathbf{K} defined by

$$(\mathbf{K}f)_p = \frac{1}{c_p} \int_0^{L_2} \Phi_p(y) f(y) dy$$

$$c_p^2 = \int_0^{L_2} \Phi_p^2(y) dy$$

then $\mathbf{K}\mathbf{A} = \mathbf{A}_p\mathbf{K}$, where

$$\mathbf{A}_p = \begin{bmatrix} 0 & 1 \\ \mu_p - j\eta & 0 \end{bmatrix}$$

From this point on, the calculations are identical to the zero-boundary condition case, except that \mathbf{K} replaces \mathbf{S} .

To consider problems where the x boundary conditions are more general than (2.2a), we assume that the matrices V_0 and V_L are now linear operators \mathbf{V}_0 and \mathbf{V}_L such that under the transform \mathbf{K} described above, the action of the operators \mathbf{V}_0 and \mathbf{V}_L are multiplicative, that is, the following transform relations hold

$$\mathbf{V}_0 m(0, y) \Leftrightarrow V_0(p) m(0, p)$$

$$\mathbf{V}_L m(L_1, y) \Leftrightarrow V_L(p) m(L_1, p)$$

where $V_0(p)$ and $V_L(p)$ are complex valued 2×2 matrix functions of p . As before, it is necessary that $V_0(p) + V_L(p) e^{A_p L_1}$ be invertible for all p . An example of a dissipative boundary condition occurs when

$$V_0(p) = \begin{bmatrix} s & 1 \\ 0 & 0 \end{bmatrix}, \quad V_L(p) = \begin{bmatrix} 0 & 0 \\ s & -1 \end{bmatrix}$$

and s is complex. This case corresponds to a damped, elastically-braced membrane on the $x = 0$ and $x = L_1$ sides. The determinant of $V_0 + V_L e^{A_p L_1}$ is

$$\left(\frac{s^2}{\sqrt{k^2 - p^2}} + \sqrt{k^2 - p^2} \right) \sin(\sqrt{k^2 - p^2} L_1) + 2s \cos(\sqrt{k^2 - p^2} L_1) \quad (5.5)$$

Typically (5.5) is non-zero.

6. Complexity Analysis

For computational purposes the rectangle Ω is discretized into an $N \times M$ grid. The complexity of the principal steps needed to calculate the least-squares estimate of the wave amplitude at one wavenumber k is given in Table 1.

TABLE 1 The Smoothing Complexity for the 2-D Helmholtz Equation

Step	Complexity	
	zero or periodic y-bound condns	other y- b.c.
Find $\Phi_p(y)$ and μ_p for (5.3),(5.4)	n/a	$O(M^2 b)$
Fourier transform the observations	$O(NM \log M)$	$O(NM^2)$
Find $\hat{m}(x, p)$ via (4.5),(4.6),(4.2)	$O(NM)$	$O(NM)$
Inverse transform $\hat{m}(x, p)$	$O(NM \log M)$	$O(NM^2)$

In Table 1, b is the bandwidth of the matrix used to approximate the operator Q in the eigenvalue calculations. Typically $b = 1$. The overall complexity for each wavenumber is then either $O(NM \log M)$ or $O(NM^2)$, depending on the y -boundary conditions. Obviously if one does not exploit the particular structure of our model, and uses only first and second moment information, the complexity would be $O(M^3 N^3)$. As discussed in the Introduction, for zero boundary conditions a com-

plexity of $O(MN \log MN)$ has been achieved in [2],[3],[4] using different techniques. The complexity of transforming the time domain images into the frequency domain for T wavenumbers is $O(TNM \log T)$. Therefore, the entire smoothing procedure would have a complexity of $O(TNM \log TM)$ or $O(TNM \log T + TNM^2)$ depending on the y-boundary conditions and assuming that $b = 1$.

7. The Smoothing Error Covariance

Using results in [1], we can express the smoothing error $\tilde{m}(x,y) = m(x,y) - \hat{m}(x,y)$ as the solution of the Hamiltonian system

$$\frac{\partial}{\partial x} \begin{bmatrix} \tilde{m}(x,y) \\ \lambda(x,y) \end{bmatrix} = \begin{bmatrix} A & qBB' \\ r^{-1}C'C & -A^* \end{bmatrix} \begin{bmatrix} \tilde{m}(x,y) \\ \lambda(x,y) \end{bmatrix} + \begin{bmatrix} B & 0 \\ 0 & C'r^{-1} \end{bmatrix} \begin{bmatrix} \epsilon(x,y) \\ w(x,y) \end{bmatrix} \quad (7.1a)$$

$$\tilde{m} \in D(A), \lambda \in D(A^*)$$

with boundary conditions (compare with (3.2))

$$V^* \Pi_v^{-1} v(y) = V^* \Pi_v^{-1} V \begin{bmatrix} \tilde{m}(0,y) \\ \tilde{m}(L_1,y) \end{bmatrix} + \begin{bmatrix} -I & 0 \\ 0 & I \end{bmatrix} \begin{bmatrix} \lambda(0,y) \\ \lambda(L_1,y) \end{bmatrix}. \quad (7.1b)$$

Transforming (7.1) with respect to y gives (see Section 3)

$$\frac{\partial}{\partial x} \begin{bmatrix} m(x,p) \\ \lambda(x,p) \end{bmatrix} = \begin{bmatrix} A_p & qBB' \\ r^{-1}C'C & -A_p^* \end{bmatrix} \begin{bmatrix} \tilde{m}(x,p) \\ \lambda(x,p) \end{bmatrix} + \begin{bmatrix} B & 0 \\ 0 & C'r^{-1} \end{bmatrix} \begin{bmatrix} \epsilon(x,p) \\ w(x,p) \end{bmatrix} \quad (7.2a)$$

$$V^* \Pi_v^{-1} v(p) = W_0 \begin{bmatrix} \tilde{m}(0,p) \\ \lambda(0,p) \end{bmatrix} + W_L \begin{bmatrix} \tilde{m}(L_1,p) \\ \lambda(L_1,p) \end{bmatrix} \quad (7.2b)$$

Our original statistical assumptions and the properties of

the transform imply that $\epsilon(x,p), \epsilon(x,q), w(x,p), w(x,q), v(p), v(q)$ are mutually uncorrelated for $p \neq q$. Therefore, if we let

$$P(x,y) = E[\bar{m}(x,y)\bar{m}^*(x,y)]$$

$$P(x,p) = E[\bar{m}(x,p)\bar{m}^*(x,p)]$$

then

$$P(x,y) = \sum_{l=-\infty}^{\infty} P(x,p) \sin^2 py, p = \frac{2\pi l}{L_2}$$

Using a Green's function solution to the acausal linear system (7.2), one can show that $P(x,p)$ satisfies

$$P(x,p) = [I|0] \Delta(x,p) \begin{bmatrix} I \\ \vdots \\ 0 \end{bmatrix}$$

where

$$\begin{aligned} \Delta(x,p) &= \int_0^{L_1} G(x,\sigma,p) \begin{bmatrix} qBB' & 0 \\ 0 & r^{-1}C'C \end{bmatrix} G^*(x,\sigma,p) d\sigma \\ &\quad + e^{H_p x} F_H^{-1} V^* \Pi_v^{-1} V (F_H^*)^{-1} e^{H_p^* x} \end{aligned} \quad (7.3)$$

and

$$G(x,\sigma,p) = \begin{cases} e^{H_p x} F_H^{-1} W_0 e^{-H_p \sigma} & x > \sigma \\ e^{H_p x} F_H^{-1} W_L e^{H_p(L_1-\sigma)} & x < \sigma \end{cases}$$

The overall complexity of solving these equations is $O(NM^2)$, due to the integration necessary to compute $\Delta(x,p)$ in (7.3). An alternate procedure with complexity $O(NM)$ can be derived by first diagonalizing (7.2) in a manner identical to that in Section 3.4. If we change variables using

$$\begin{bmatrix} e_f(x, p) \\ e_b(x, p) \end{bmatrix} = M_p^{-1} \begin{bmatrix} \tilde{m}(x, p) \\ -\lambda(x, p) \end{bmatrix}$$

then

$$\frac{\partial}{\partial x} \begin{bmatrix} e_f(x, p) \\ e_b(x, p) \end{bmatrix} = \begin{bmatrix} \Lambda_p & 0 \\ 0 & -\Lambda_p \end{bmatrix} \begin{bmatrix} e_f(x, p) \\ e_b(x, p) \end{bmatrix} + \begin{bmatrix} B_f^\epsilon \\ B_b^\epsilon \end{bmatrix} \begin{bmatrix} \epsilon(x, p) \\ w(x, p) \end{bmatrix}$$

$$V^* \Pi_v^{-1} v(p) = \begin{bmatrix} V_f^0 & V_b^0 \\ e_b(0, p) \end{bmatrix} \begin{bmatrix} e_f(0, p) \\ e_b(0, p) \end{bmatrix} + \begin{bmatrix} V_f^L & V_b^L \\ e_b(L_1, p) \end{bmatrix} \begin{bmatrix} e_f(L_1, p) \\ e_b(L_1, p) \end{bmatrix}$$

where

$$\begin{bmatrix} B_f^\epsilon \\ B_b^\epsilon \end{bmatrix} = M_p^{-1} \begin{bmatrix} B & 0 \\ 0 & C' r^{-1} \end{bmatrix}$$

A solution to these equations is

$$\begin{bmatrix} e_f(x, p) \\ e_b(x, p) \end{bmatrix} = \begin{bmatrix} e^{\Lambda_p x} & 0 \\ 0 & e^{\Lambda_p (L_1 - x)} \end{bmatrix} F_{fb}^{-1} \{ V^* \Pi_v^{-1} v(p) - V_f^L e_f^0(L_1, p) - V_b^0 e_b^0(0, p) \} \\ + \begin{bmatrix} e_f^0(x, p) \\ e_b^0(x, p) \end{bmatrix} \quad (7.4)$$

where

$$\frac{\partial}{\partial x} e_f^0(x, p) = \Lambda_p e_f^0(x, p) + B_f^\epsilon \begin{bmatrix} \epsilon(x, p) \\ w(x, p) \end{bmatrix} \quad (7.5a)$$

$$\frac{\partial}{\partial x} e_b^0(x, p) = -\Lambda_p e_b^0(x, p) + B_b^\epsilon \begin{bmatrix} \epsilon(x, p) \\ w(x, p) \end{bmatrix} \quad (7.5b)$$

$$e_f^0(0, p) = 0, e_b^0(L_1, p) = 0 \quad (7.5c)$$

We can now express $\Delta(x, p)$ as

$$\Delta(x, p) = M_p \Theta(x, p) M_p^* \quad (7.6)$$

where

$$\Theta(x, p) = E \begin{bmatrix} e_f(x, p) \\ e_b(x, p) \end{bmatrix} [e_f^*(x, p) \quad e_b^*(x, p)]$$

We will express the diagonalized error covariance $\Theta(x, p)$ in terms of the second moments of the random variables $\{v, e_f^0(x), e_b^0(x)\}$, where we have suppressed the argument p for simplicity. Define the covariances

$$R_f(x_1, x_2) = E[e_f^0(x_1)e_f^{0*}(x_2)]$$

$$R_b(x_1, x_2) = E[e_b^0(x_1)e_b^{0*}(x_2)]$$

$$R_{fb}(x_1, x_2) = E[e_f^0(x_1)e_b^{0*}(x_2)]$$

Now $\Theta(x, p)$ can be written as

$$\Phi_{fb}(x) F_{fb}^{-1} \left\{ \frac{1}{2L_2} V^* \Pi_v^{-1} V + V_f^L \Pi_f(L_1) V_f^{L*} + V_f^L R_{fb}(L_1, 0) V_b^{0*} + \right.$$

$$V_b^0 \Pi_b(0) V_b^{0*} + V_b^0 R_{fb}^*(L_1, 0) V_f^{L*} \right\} (F_{fb}^*)^{-1} \Phi_{fb}^*(x)$$

$$- \Phi_{fb}(x) F_{fb}^{-1} G(x) - G^*(x) (F_{fb}^*)^{-1} \Phi_{fb}^*(x) + \begin{bmatrix} \Pi_f(x) & 0 \\ 0 & \Pi_b(x) \end{bmatrix}$$

where

$$\Phi_{fb}(x) = \begin{bmatrix} e^{\Lambda_p x} & 0 \\ 0 & e^{\Lambda_p(L_1-x)} \end{bmatrix}$$

$$G(x) = V_f^L \left[R_f(L_1, x) : R_{fb}(L_1, x) \right] + V_b^0 \left[R_{fb}^*(x, 0) : R_b(0, x) \right]$$

$$R_f(L_1, x) = e^{\Lambda_p(L_1-x)} \Pi_f(x) , \quad R_b(0, x) = e^{\Lambda_p x} \Pi_b(x)$$

$$R_{fb}(x_1, x_2) = -e^{\Lambda_p(x_1-x_2)} \Pi_{fb}(x_1-x_2) , \quad x_1 > x_2$$

with

$$\frac{\partial}{\partial x} \Pi_f(x) = \Lambda_p \Pi_f(x) + \Pi_f(x) \Lambda_p^* + \frac{1}{2L_2} B_f^q \begin{bmatrix} q & 0 \\ 0 & r \end{bmatrix} B_f^r , \quad \Pi_f(0) = 0$$

$$- \frac{\partial}{\partial x} \Pi_b(x) = \Lambda_p \Pi_b(x) + \Pi_b(x) \Lambda_p^* + \frac{1}{2L_2} B_b^q \begin{bmatrix} q & 0 \\ 0 & r \end{bmatrix} B_b^r , \quad \Pi_b(L_1) = 0$$

$$\Pi_{fb}(x) = \frac{1}{2L_2} \int_0^x e^{-\Lambda_f u} B_f^e \begin{bmatrix} q & 0 \\ 0 & r \end{bmatrix} B_f^e * e^{\Lambda_f u} du$$

8. Concluding Remarks

In this paper, we have constructed a recursive smoother for the 2-D Helmholtz equation. The advantage of our method is that it leads to computationally attractive algorithms even when the x-boundary conditions are dissipative and include random inputs. Another advantage is that a linear least-squares Born inversion of the wavefield can be performed along with the image restoration without changing the computational complexity. Higher order equations characterizing vibrating plates, etc. could be approached in a manner similar to that presented here. If the wavenumber k has variations in the x-direction, then an approach based on operator Riccati equations could be used to formally diagonalize the smoother dynamics. However, the initial value problem

$$\frac{\partial}{\partial x} m(x, y) = A m(x, y)$$

$$m(0, y) = m_0(y)$$

is not well-posed (it does not lead to a semigroup), and this raises questions about the existence and uniqueness of solutions to the corresponding Riccati equations.

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Appendix A

In this appendix, we show that the state variable representation (3.1) of the Helmholtz equation is well-posed. To do this, we show that the transformed system (3.4) is well-posed, and that the solutions to these acausal linear systems give rise to Fourier coefficients that ensure that the formal Fourier series does converge and that the state vector $m(x,y) \in D(A)$.

(3.4) are well posed for every p because the matrix

$$V_0 + V_L e^{A, L_1} = \begin{bmatrix} 1 & 0 \\ \cosh \beta L_1 & \beta^{-1} \sinh(\beta L_1) \end{bmatrix} \quad (\text{A.1})$$

where $\beta = (p^2 - k^2)^{1/2}$, is invertible for all p . If we now change variables in (3.4) as follows:

$$\begin{bmatrix} \Upsilon_f(x,p) \\ \Upsilon_b(x,p) \end{bmatrix} = D_p^{-1} \begin{bmatrix} m_1(x,p) \\ m_2(x,p) \end{bmatrix} \quad (\text{A.2})$$

$$D_p = \begin{bmatrix} 1 & 1 \\ -\beta & \beta \end{bmatrix}$$

then

$$\frac{\partial}{\partial x} \begin{bmatrix} \Upsilon_f(x,p) \\ \Upsilon_b(x,p) \end{bmatrix} = \begin{bmatrix} -\beta & 0 \\ 0 & \beta \end{bmatrix} \begin{bmatrix} \Upsilon_f(x,p) \\ \Upsilon_b(x,p) \end{bmatrix} + 1/2 \begin{bmatrix} -\beta^{-1} \\ \beta^{-1} \end{bmatrix} \epsilon(x,p) \quad (\text{A.3a})$$

$$v(p) = \begin{bmatrix} 1 & 1 \\ 0 & 0 \end{bmatrix} \begin{bmatrix} \Upsilon_f(0, p) \\ \Upsilon_b(0, p) \end{bmatrix} + \begin{bmatrix} 0 & 0 \\ 1 & 1 \end{bmatrix} \begin{bmatrix} \Upsilon_f(L_1, p) \\ \Upsilon_b(L_1, p) \end{bmatrix} \quad (\text{A.3b})$$

We can write the solution to (A.3) as

$$\begin{bmatrix} \Upsilon_f(x, p) \\ \Upsilon_b(x, p) \end{bmatrix} = (1 - e^{-2\beta L_1})^{-1} \begin{bmatrix} e^{-\beta x} & 0 \\ 0 & e^{-\beta(L_1-x)} \end{bmatrix} \begin{bmatrix} 1 & -e^{-\beta L_1} \\ -e^{-\beta L_1} & 1 \end{bmatrix} \begin{bmatrix} v_1(p) - \Upsilon_b^0(0, p) \\ v_2(p) - \Upsilon_f^0(L_1, p) \end{bmatrix} \\ + \begin{bmatrix} \Upsilon_f^0(x, p) \\ \Upsilon_b^0(x, p) \end{bmatrix} \quad (\text{A.4})$$

where

$$\frac{\partial}{\partial x} \Upsilon_f^0(x, p) = -\beta \Upsilon_f^0(x, p) - 1/2 \beta^{-1} \epsilon(x, p)$$

$$\frac{\partial}{\partial x} \Upsilon_b^0(x, p) = \beta \Upsilon_b^0(x, p) + 1/2 \beta^{-1} \epsilon(x, p)$$

$$\Upsilon_f^0(0, p) = \Upsilon_b^0(L_1, p) = 0$$

$\Upsilon_f(x, p)$ and $\Upsilon_b(x, p)$ are continuous functions of x that depend continuously on v and ϵ . Equations (3.1) will have a unique solution which depends continuously on v and ϵ if the appropriate Fourier series converge and if $m(x, y) \in D(A)$. We will assume that $(\partial^2/\partial^2 y)v(y) \in \mathbf{L}_2^2(0, L_1)$, $\epsilon(x, y) \in \mathbf{L}_2(\Omega)$ and show that this implies that $m_1(x, y) \in D(T)$ and $m_2(x, y) \in \mathbf{L}_2(\Omega)$. We first derive bounds for $\Upsilon_f^0(x, p)$ and $\Upsilon_b^0(x, p)$. Solving for $\Upsilon_f^0(x, p)$ gives

$$\Upsilon_f^0(x, p) = \frac{-1}{2\beta} \int_0^x e^{-\beta(x-u)} \epsilon(u, p) du$$

which gives

$$|\Upsilon_f^0(x, p)| \leq \frac{1}{2|\beta|} \left[\int_0^x |e^{-\beta(x-u)}|^2 du \right]^{1/2} \left[\int_0^x |\epsilon(u, p)|^2 du \right]^{1/2}$$

where

$$\int_0^x |e^{-\beta(x-y)}|^2 dy \leq \frac{1}{2R\epsilon(\beta)}$$

Thus

$$|\Upsilon_f^0(x, p)| \leq \frac{1}{2|\beta|\sqrt{2R\epsilon(\beta)}} \left[\int_0^x |\epsilon(u, p)|^2 du \right]^{1/2}$$

Since $\beta \approx p$ when $p \rightarrow \infty$

$$|\Upsilon_f^0(x, p)|^2 \sim O(Gp^{-3})$$

where

$$G = \int_0^{L_1} |\epsilon(u, p)|^2 du$$

In a similar manner,

$$|\Upsilon_b^0(x, p)|^2 \sim O(Gp^{-3})$$

Our assumption on the input field ϵ implies that

$$\int_0^{L_2} |\epsilon(x, y)|^2 dy < \infty$$

almost everywhere with respect to x (a.e. wrt x). By Parseval's theorem

$$\sum_p |\epsilon(x, p)|^2 < \infty \text{ a.e. wrt } x$$

and hence by the comparison test for series

$$|\epsilon(x, p)|^2 \sim o(p^{-1})$$

a.e. wrt x . Therefore

$$|\Upsilon_f^0(x, p)|^2 \sim |\Upsilon_b^0(x, p)|^2 \sim O(p^{-4})$$

The first term in (A.4) involving the boundary term $v(p)$ is becoming exponentially small as $p \rightarrow \infty$ for all x in $(0, L_1)$. Since $(\partial^2/\partial^2 y)v(y) \in L_2^2(0, L_1)$ by assumption, the first term in (A.4) has a convergent Fourier series that is

twice differentiable with respect to y on $[0, L_1]$, and is infinitely differentiable in $(0, L_1)$. The second term of (A.4) therefore dominates the solution in $(0, L_1)$ as p gets large, hence

$$|\Upsilon_f(x, p)|^2 \sim |\Upsilon_b(x, p)|^2 \sim O(p^{-4})$$

a.e. wrt x . By using the change of variables (A.1), we see that

$$|m_1(x, p)|^2 = |\Upsilon_f(x, p) + \Upsilon_b(x, p)|^2 \sim O(p^{-4})$$

and

$$|m_2(x, p)|^2 = |-\beta\Upsilon_f(x, p) + \beta\Upsilon_b(x, p)|^2 \sim O(p^{-2})$$

These bounds show that
 $m_1(x, y) \in D(T)$ and $m_2(x, y) \in L_2(\Omega)$ [5].

Appendix B

We show in this appendix that for the 2-D Helmholtz equation, when the x -boundary conditions are separable ($\theta_{12} = \theta_{21} = 0$), and the observed image is bounded and mean-square differentiable in the y -direction,

$$\hat{m}(x, y) \in D(A) \text{ and } \hat{\lambda}(x, y) \in D(A^*)$$

To begin with, it is sufficient to show that $(\partial^2/\partial^2 y)\Psi_f(x, y)$ and $(\partial^2/\partial^2 y)\Psi_b(x, y)$ exist in the mean-square sense. The boundary conditions will be satisfied due to the sine transform.

The p -domain equations for $\Psi_{f1}^0, \Psi_{f2}^0, \Psi_{b1}^0, \Psi_{b2}^0$ are

given by ($\Psi_{f,i}(x,y)$ denotes the i th component of $\Psi_f(x,y)$, etc.)

$$\Psi_{f1}^0(x,p) = \int_0^x e^{\lambda_0(x-s)} j \frac{q}{4\lambda_0 r \sigma} z(s,p) ds$$

$$\Psi_{f2}^0(x,p) = - \int_0^x e^{\bar{\lambda}_0(x-s)} j \frac{q}{4\bar{\lambda}_0 r \sigma} z(s,p) ds$$

$$\Psi_{b1}^0(x,p) = \int_x^{L_1} e^{-\lambda_0(x-s)} j \frac{q}{4\lambda_0 r \sigma} z(s,p) ds$$

$$\Psi_{b2}^0(x,p) = - \int_x^{L_1} e^{-\bar{\lambda}_0(x-s)} j \frac{q}{4\bar{\lambda}_0 r \sigma} z(s,p) ds$$

As in Appendix A, one can show that

$$\|\Psi_b^0(x,p)\|_{\mathbf{R}^2} < g_1 p^{-3/2} M(p) \quad (\text{B.1})$$

where

$$M^2(p) = \int_0^{L_1} |z(x,p)|^2 dx$$

and g_1 is a finite constant. The same result holds for Ψ_f^0 .

For separable boundary conditions [1]

$$\Psi_f(x,p) = \Psi_f^B(x,p) + \Psi_f^0(x,p)$$

Where for large p

$$\Psi_f^B(x,p) = \left[e^{\Lambda_p x} |0 \right] F_B^{-1} V_B^0 \Psi_B^0(0,p)$$

The $L^2(\Omega)$ norm of $\Psi_f(x,y)$ can be expressed using Parseval's relation as

$$\|\Psi_f(x,y)\|_{L^2(\Omega)}^2 = \sum_{p=1}^{\infty} \int_0^{L_1} \|\Psi_f(x,p)\|_{\mathbf{R}^2}^2 dx$$

In a similar fashion:

$$\left\| \frac{\partial^2}{\partial^2 y} \Psi_f(x, y) \right\|_{L^2(\Omega)}^2 = \sum_{p=1}^{\infty} p^4 \int_0^{L_1} \|\Psi_f(x, p)\|_{R^2}^2 dx \quad (\text{B.2})$$

We wish to prove that this norm is finite. For separable boundary conditions one can verify that

$$\|F_{fb}^{-1} V_b^0\| < g_2 p$$

and

$$\left\| e^{\Lambda_0 x} z(0) \right\| < g_3 e^{\operatorname{Re} \lambda_0 x}$$

Combining these bounds gives

$$\|\Psi_f^B(x, p)\|_{R^2} < g_4 p^{-1/2} M(p) e^{\operatorname{Re} \lambda_0 x}$$

so that

$$\left\| \frac{\partial^2}{\partial^2 y} \Psi_f^B(x, y) \right\|_{L^2(\Omega)}^2 < \sum_{p=1}^{\infty} p^3 g_4 M^2(p) \int_0^{L_1} e^{2\operatorname{Re} \lambda_0 x} dx$$

or

$$\left\| \frac{\partial^2}{\partial^2 y} \Psi_f^B(x, y) \right\|_{L^2(\Omega)}^2 < \sum_{p=1}^{\infty} p^2 g_5 M^2(p)$$

since $\lambda_0 < g_6 p$. If we assume that $(\partial/\partial y)z(x, y) \in L^2(\Omega)$ then

$$\left\| \frac{\partial}{\partial y} z(x, y) \right\|_{L^2(\Omega)} = \sum_{p=1}^{\infty} p^2 M^2(p) < \infty$$

which shows that $\Psi_f^B(x, y)$ is twice differentiable in the y -direction. Similarly,

$$\left\| \frac{\partial^2}{\partial^2 y} \Psi_f^0(x, y) \right\|_{L^2(\Omega)} < \sum_{p=1}^{\infty} L_1 p M^2(p) < \infty$$

which implies that $\Psi_f(x, y)$ has a mean-square second derivative in the y -direction. A similar argument shows that $\Psi_b(x, y)$ is also mean-square twice differentiable in the y -direction.

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8a. NAME OF FUNDING/SPONSORING ORGANIZATION	8b. OFFICE SYMBOL <i>(If applicable)</i>	9. PROCUREMENT INSTRUMENT IDENTIFICATION NUMBER <i>N00014-85-K-0255</i>										
8c. ADDRESS (City, State and ZIP Code)		10. SOURCE OF FUNDING NOS. <table border="1"><tr><td>PROGRAM ELEMENT NO.</td><td>PROJECT NO.</td><td>TASK NO.</td><td>WORK UNIT NO.</td></tr><tr><td></td><td></td><td></td><td><i>NR661-019</i></td></tr></table>		PROGRAM ELEMENT NO.	PROJECT NO.	TASK NO.	WORK UNIT NO.				<i>NR661-019</i>	
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			<i>NR661-019</i>									
11. TITLE (Include Security Classification) <i>Recursive Linear Smoothing for the 2-D Helmholtz Equation (Unclassified)</i>												
12. PERSONAL AUTHOR(S) <i>Riddle, L. R. and Weinert, H. L.</i>												
13a. TYPE OF REPORT <i>Interim</i>	13b. TIME COVERED <i>FROM 5/1/85 TO 9/22/86</i>	14. DATE OF REPORT (Yr., Mo., Day) <i>September 25, 1986</i>	15. PAGE COUNT <i>28</i>									
16. SUPPLEMENTARY NOTATION												
17. COSATI CODES <table border="1"><tr><td>FIELD</td><td>GROUP</td><td>SUB. GR.</td></tr><tr><td></td><td></td><td></td></tr><tr><td></td><td></td><td></td></tr></table>		FIELD	GROUP	SUB. GR.							18. SUBJECT TERMS (Continue on reverse if necessary and identify by block number) <i>Helmholtz equation, wave equation, smoothing, image processing, distributed parameter systems, acausal systems.</i>	
FIELD	GROUP	SUB. GR.										
19. ABSTRACT (Continue on reverse if necessary and identify by block number) <i>A fast algorithm for reconstructing images governed by a 2-D Helmholtz equation is presented. The computational complexity is $O(NM\log M)$ or $O(NM^2)$, depending on boundary conditions, where N and M are the number of spatial grid points in the x and y directions, respectively. We show that the smoothing problem is well-posed, and that the sample functions of the smoothed estimate possess smoothness properties consistent with the Helmholtz equation.</i>												
20. DISTRIBUTION/AVAILABILITY OF ABSTRACT <input checked="" type="checkbox"/> UNCLASSIFIED/UNLIMITED <input type="checkbox"/> SAME AS RPT <input type="checkbox"/> DTIC USERS <input type="checkbox"/>		21. ABSTRACT SECURITY CLASSIFICATION <i>Unclassified</i>										
22a. NAME OF RESPONSIBLE INDIVIDUAL <i>Dr. Neil L. Gerr</i>		22b. TELEPHONE NUMBER <i>(Include Area Code)</i> <i>(202) 696-4321</i>	22c. OFFICE SYMBOL									

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